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# Linear rigidity of stationary stochastic processes

Alexander I. Bufetov <sup>\*</sup>    Yoann Dabrowski <sup>†</sup>    Yanqi Qiu <sup>‡</sup>

## Abstract

We consider stationary stochastic processes  $X_n$ ,  $n \in \mathbb{Z}$  such that  $X_0$  lies in the closed linear span of  $X_n$ ,  $n \neq 0$ ; following Ghosh and Peres, we call such processes linearly rigid. Using a criterion of Kolmogorov, we show that it suffices, for a stationary stochastic process to be rigid, that the spectral density vanish at zero and belong to the Zygmund class  $\Lambda_*(1)$ . We next give sufficient condition for stationary determinantal point processes on  $\mathbb{Z}$  and on  $\mathbb{R}$  to be rigid. Finally, we show that the determinantal point process on  $\mathbb{R}^2$  induced by a tensor square of Dyson sine-kernels is *not* linearly rigid.

## 1 Introduction

This paper is devoted to rigidity of stationary determinantal point processes.

Recall that stationary determinantal point processes are strongly chaotic: they have the Kolmogorov property (Lyons [9]) and the Bernoulli property (Lyons and Steif [10]); and they satisfy the Central Limit Theorem (Costin and Lebowitz [2], Soshnikov [13]). On the other hand, Ghosh [5] and Ghosh-Peres [6] proved, for the determinantal point processes such as Dyson sine process and Ginibre point process, that number of particles in a finite window is measurable with respect to the

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completion of the sigma-algebra describing the configurations outside that finite window. Their argument is spectral: they construct, for any small  $\varepsilon$ , a compactly supported smooth function  $\varphi_\varepsilon$ , such that  $\varphi_\varepsilon$  equals 1 in a fixed finite window and the linear statistic corresponding to  $\varphi_\varepsilon$  has variance smaller than  $\varepsilon$ .

In the same spirit, we consider general stationary stochastic processes (in broad sense)  $X_n, n \in \mathbb{Z}$  such that  $X_0$  lies in the closed linear span of  $X_n, n \neq 0$ ; following Ghosh and Peres, we call such processes linearly rigid. Using a criterion of Kolmogorov, we show that it suffices, for a stationary stochastic process to be rigid, that the spectral density vanish at zero and belong to the Zygmund class  $\Lambda_*(1)$ . We next give sufficient condition for stationary determinantal point processes on  $\mathbb{Z}$  and on  $\mathbb{R}$  to be rigid. Finally, we show that the determinantal point process on  $\mathbb{R}^2$  induced by a tensor square of Dyson sine-kernels is *not* linearly rigid.

We now turn to more precise statements. Let  $X = \{X_n : n \in \mathbb{Z}^d\}$  be a multi-dimensional time stationary stochastic process of real-valued random variables defined on a probability space  $(\Omega, \mathbb{P})$ . Let  $H(X) \subset L^2(\Omega, \mathbb{P})$  denote the closed subspace linearly spanned by  $\{X_n : n \in \mathbb{Z}^d\}$  and let  $\check{H}_0(X)$  denote the one linearly spanned by  $\{X_n : n \in \mathbb{Z}^d \setminus \{0\}\}$ .

**Definition 1.1.** The stochastic process  $X$  is said to be linearly rigid if

$$X_0 \in \check{H}_0(X). \quad (1)$$

Let  $\text{Conf}(\mathbb{R}^d)$  be the set of configurations on  $\mathbb{R}^d$ . For a bounded Borel subset  $B \subset \mathbb{R}^d$ , we denote  $N_B : \text{Conf}(\mathbb{R}^d) \rightarrow \mathbb{N} \cup \{0\}$  the function defined by

$$N_B(\mathcal{X}) := \text{the cardinality of } B \cap \mathcal{X}.$$

The space  $\text{Conf}(\mathbb{R}^d)$  is equipped with the Borel  $\sigma$ -algebra which is the smallest  $\sigma$ -algebra making all  $N_B$ 's measurable. Recall that a point process with phase space  $\mathbb{R}^d$  is, by definition, a Borel probability measure on the space  $\text{Conf}(\mathbb{R}^d)$ . For the background on point process, the reader is referred to Daley and Vere-Jones' book [3].

Given a stationary point process on  $\mathbb{R}^d$  and  $\lambda > 0$ , we introduce the stationary stochastic process  $N^{(\lambda)} = (N_n^{(\lambda)})_{n \in \mathbb{Z}^d}$  by the formula

$$N_n^{(\lambda)}(\mathcal{X}) := \text{the cardinality of } \mathcal{X} \cap (n\lambda + [-\lambda/2, \lambda/2)^d). \quad (2)$$

**Definition 1.2.** A stationary point process  $\mathbb{P}$  on  $\mathbb{R}^d$  is called **linearly rigid**, if for any  $\lambda > 0$ , the stationary stochastic process  $N^{(\lambda)} = (N_n^{(\lambda)})_{n \in \mathbb{Z}^d}$  is linearly rigid, i.e.,

$$N_0^{(\lambda)} \in \check{H}_0(N^{(\lambda)}).$$

The above definition is motivated by the definition due to Ghosh and Peres of rigidity of point processes on  $\mathbb{R}^d$ , see [5] and [6]. Given a Borel subset  $C \subset \mathbb{R}^d$ , we will denote

$$\mathcal{F}_C = \sigma(\{N_B : B \subset C, B \text{ bounded Borel}\})$$

the  $\sigma$ -algebra generated by all random variables of the form  $N_B$  where  $B \subset C$  ranges over all bounded Borel subsets of  $C$ . Let  $\mathbb{P}$  be a point process on  $\mathbb{R}^d$ , i.e.,  $\mathbb{P}$  is a Borel probability on  $\text{Conf}(\mathbb{R}^d)$ , and denote  $\mathcal{F}_C^\mathbb{P}$  for the  $\mathbb{P}$ -completion of  $\mathcal{F}_C$ .

**Definition 1.3** (Ghosh [5], Ghosh-Peres [6]). A point process  $\mathbb{P}$  on  $\mathbb{R}^d$  is called **rigid**, if for any bounded Borel set  $B \subset \mathbb{R}^d$  with Lebesgue-negligible boundary  $\partial B$ , the random variable  $N_B$  is  $\mathcal{F}_{\mathbb{R}^d \setminus B}^\mathbb{P}$ -measurable.

*Remark 1.1.* Of course, in the above definition, it suffices to take Borel sets  $B$  of the form  $[-\gamma, \gamma]^d$  for  $\gamma > 0$ , cf. [6].

A linear rigid stationary point process on  $\mathbb{R}^d$  is of course rigid in the sense of Ghosh and Peres. Observe that proofs for rigidity in [5], [6] and [1] in fact establish linear rigidity. We would like also to mention a notion of insertion-deletion tolerance studied by Holroyd and Soo in [7], which is in contrast to the notion of rigidity property.

## 2 The Kolmogorov criterion for linear rigidity

In this note, the Fourier transform of a function  $f : \mathbb{R}^d \rightarrow \mathbb{C}$  is defined as

$$\hat{f}(\xi) = \int_{\mathbb{R}^d} f(x) e^{-i2\pi x \cdot \xi} dx.$$

Denote by  $\mathbb{T}^d = \mathbb{R}^d / \mathbb{Z}^d$  the  $d$ -dimensional torus. In what follows, we identify  $\mathbb{T}^d$  with  $[-1/2, 1/2]^d$ . The Fourier coefficients of a measure  $\mu$  on  $\mathbb{T}^d$  are given, for any  $k \in \mathbb{Z}^d$ , by the formula

$$\hat{\mu}(k) = \int_{\mathbb{T}^d} e^{-i2\pi k \cdot \theta} d\mu_X(\theta), \text{ where } k \cdot \theta := k_1\theta_1 + \dots + k_d\theta_d.$$

Denote by  $\mu_X$  the spectral measure of  $X$ , i.e.,

$$\forall k \in \mathbb{Z}^d, \quad \mathbb{E}(X_0 X_k) = \mathbb{E}(X_n X_{n+k}) = \int_{\mathbb{T}^d} e^{-i2\pi k \cdot \theta} d\mu_X(\theta) = \hat{\mu}_X(k). \quad (3)$$

Recall that we have the following natural isometric isomorphism

$$H(X) \simeq L^2(\mathbb{T}^d, \mu_X), \quad (4)$$

by assigning to  $X_n \in H(X)$  the function  $\theta \mapsto e^{i2\pi n \cdot \theta} \in L^2(\mathbb{T}^d, \mu_X)$ .

Let  $\mu_X = \mu_a + \mu_s$  be the Lebesgue decomposition of  $\mu_X$  with respect to the normalized Lebesgue measure  $m(d\theta) = d\theta_1 \cdots d\theta_d$  on  $\mathbb{T}^d$ , i.e.,  $\mu_a$  is absolutely continuous with respect to  $m$  and  $\mu_s$  is singular to  $m$ . Set

$$\omega_X(\theta) := \frac{d\mu_a}{dm}(\theta).$$

**Lemma 2.1** (The Kolmogorov Criterion). *We have*

$$\text{dist}(X_0, \check{H}_0(X)) = \left( \int_{\mathbb{T}^d} \omega_X^{-1} dm \right)^{-1/2}.$$

*The right side is to be interpreted as zero if  $\int_{\mathbb{T}^d} \omega_X^{-1} dm = \infty$ .*

When the measure  $\mu$  is assumed to be absolutely continuous with respect to  $m$ , Lemma 2.1 is a result of Kolmogorov, see Remark 5.17 in Lyons-Steif [10].

**Corollary 2.2.** *The stationary stochastic process  $X = (X_n)_{n \in \mathbb{Z}^d}$  is linearly rigid if and only if*

$$\int_{\mathbb{T}^d} \omega_X^{-1} dm = \infty.$$

*Proof of Lemma 2.1.* We follow the argument of Lyons-Steif [10]. By the Lebesgue decomposition of  $\mu$ , we may take a subset  $A \subset \mathbb{T}^d$  of full Lebesgue measure  $m(A) = 1$ , such that  $\mu_a(A) = 1$  and  $\mu_s(A) = 0$ .

Denote

$$L_0 = \overline{\text{span}}^{L^2(\mathbb{T}^d, \mu_X)}[e^{i2\pi n \cdot \theta} : n \neq 0].$$

By the isometric isomorphism (4), it suffices to show that

$$\text{dist}(1, L_0) = \left( \int_{\mathbb{T}^d} \omega_X^{-1} dm \right)^{-1/2}, \quad (5)$$

where 1 is the constant function taking value 1. Write

$$1 = p + h, \text{ such that } p \perp L_0, h \in L_0.$$

Modifying, if necessary, the values of  $p$  and  $h$  on a  $\mu$ -negligible subset, we may assume that

$$1 = p(\theta) + h(\theta) \text{ for all } \theta \in \mathbb{T}^d.$$

Since  $p \perp L_0$ , we have

$$0 = \langle p, e^{i2\pi n \cdot \theta} \rangle_{L^2(d\mu)} = \int_{\mathbb{T}^d} p(\theta) e^{-i2\pi n \cdot \theta} d\mu(\theta), \text{ for any } n \in \mathbb{Z}^d \setminus 0. \quad (6)$$

Therefore, the complex measure  $p \cdot d\mu$  is a multiple of Lebesgue measure, i.e., there exists  $\xi \in \mathbb{C}$ , such that

$$p \cdot d\mu = \xi dm.$$

It follows that  $p$  must vanish almost everywhere with respect to the singular component  $\mu_s$  of  $\mu$ , and  $p(\theta)\omega_X(\theta) = \xi$  for  $m$ -almost every  $\theta \in \mathbb{T}^d$ . Thus we have

$$\|p\|_{L^2(d\mu)} = \|p\|_{L^2(d\mu_a)}, \quad (7)$$

and

$$h(\theta) = 1 - \xi\omega_X(\theta)^{-1} \text{ for } m\text{-almost every } \theta \in \mathbb{T}^d. \quad (8)$$

**Case 1:**  $\int_{\mathbb{T}^d} \omega_X^{-1} dm < \infty$ .

Define a function  $f : \mathbb{T}^d \rightarrow \mathbb{C}$  by  $f = \omega_X^{-1} \chi_A$ . Then  $f \in L^2(d\mu) \ominus L_0$ . Indeed,

$$\|f\|_{L^2(d\mu)}^2 = \int_{\mathbb{T}^d} \omega_X^{-2} \chi_A d\mu = \int_{\mathbb{T}^d} \omega_X^{-2} d\mu_a = \int_{\mathbb{T}^d} \omega_X^{-1} dm < \infty.$$

And, for all  $n \in \mathbb{Z}^d \setminus 0$ ,

$$\langle f, e^{i2\pi n \cdot \theta} \rangle_{L^2(d\mu)} = \int_{\mathbb{T}^d} \omega_X(\theta)^{-1} \chi_A(\theta) e^{-i2\pi n \cdot \theta} d\mu(\theta) = \int_{\mathbb{T}^d} e^{-i2\pi n \cdot \theta} dm(\theta) = 0.$$

It follows that  $f \perp h$ , i.e.,

$$0 = \langle h, f \rangle_{L^2(d\mu)} = \int_{\mathbb{T}^d} h \omega_X^{-1} \chi_A d\mu = \int_{\mathbb{T}^d} h dm.$$

By (8), we get

$$\int_{\mathbb{T}^d} (1 - \xi \omega_X^{-1}) dm = 0,$$

and hence

$$\xi = \left( \int_{\mathbb{T}^d} \omega_X^{-1} dm \right)^{-1}.$$

It follows that

$$\text{dist}(1, L_0)^2 = \|p\|_{L^2(d\mu)}^2 = \|p\|_{L^2(d\mu_a)}^2 = \xi^2 \int_{\mathbb{T}^d} \omega_X^{-2} \omega_X dm = \xi.$$

This shows the desired equality (5).

**Case 2:**  $\int_{\mathbb{T}^d} \omega_X^{-1} dm = \infty$ .

We claim that  $\xi = 0$ . If the claim were verified, then we would get the desired identity in this case

$$\text{dist}(1, L_0) = 0.$$

So let us turn to the proof of the claim. We argue by contradiction. If  $\xi \neq 0$ , then  $p \neq 0$  and

$$\|p\|_{L^2(d\mu)}^2 = \|p\|_{L^2(d\mu_a)}^2 = \xi^2 \|\omega_X^{-1}\|_{L^2(d\mu_a)}^2 = \xi^2 \int_{\mathbb{T}^d} \omega_X^{-1} dm = \infty.$$

This contradicts the fact that  $p \in L^2(d\mu)$ . □

*Remark 2.1.* The same argument shows that, in the case of one-dimensional time, the following assertions are equivalent:

- $\sum_{k=-n}^n X_k \in \overline{\text{span}}\{X_j : |j| \geq n+1\}$ ;
- for any  $\alpha_1, \dots, \alpha_n \in (-1/2, 1/2) \setminus \{0\}$ , we have

$$\int_{\mathbb{T}} \frac{\prod_{j=1}^n |e^{i2\pi\theta} - e^{i2\pi\alpha_j}|^2 |e^{i2\pi\theta} - e^{-i2\pi\alpha_j}|^2}{\omega_X(\theta)} dm(\theta) = \infty.$$

It would be interesting to find a similar characterization for multi-dimensional time as well.

Denote by  $\text{Cov}(U, V)$  the covariance between two random variables  $U$  and  $V$ :  
 $\text{Cov}(U, V) = \mathbb{E}(UV) - \mathbb{E}(U)\mathbb{E}(V)$ .

If  $X = (X_n)_{n \in \mathbb{Z}^d}$  is a stochastic process such that

$$\sum_{n \in \mathbb{Z}^d} |\text{Cov}(X_0, X_n)| < \infty, \quad (9)$$

then we may define a continuous function on  $\mathbb{T}^d$  by the formula

$$\omega_X(\theta) := \sum_{n \in \mathbb{Z}^d} \text{Cov}(X_0, X_n) e^{i2\pi n \cdot \theta}. \quad (10)$$

**Lemma 2.3.** *Let  $X = (X_n)_{n \in \mathbb{Z}^d}$  be a stationary stochastic process satisfying condition (9). Then we have the following explicit Lebesgue decomposition of  $\mu_X$ :*

$$\mu_X = (\mathbb{E}X_0)^2 \cdot \delta_0 + \omega_X \cdot m, \quad (11)$$

where  $\delta_0$  is the Dirac measure on the point  $0 \in \mathbb{T}^d$  and  $\omega_X$  is the function on  $\mathbb{T}^d$  defined by (10).

*Proof.* Note that, under the assumption (9), the function  $\omega_X(\theta)$  is well-defined and continuous on  $\mathbb{T}^d$ . For proving the decomposition (11), it suffices to show that the Fourier coefficients of  $\mu_X$  coincide with those of  $\nu_X := (\mathbb{E}X_0)^2 \cdot \delta_0 + \omega_X \cdot m$ . But if  $n \in \mathbb{Z}^d$ , then

$$\hat{\nu}_X(n) = (\mathbb{E}X_0)^2 + \text{Cov}(X_0, X_n) = \mathbb{E}(X_0 X_n) = \hat{\mu}_X(n).$$

The lemma is completely proved.  $\square$

### 3 A sufficient condition for linear rigidity

**Theorem 3.1.** *Let  $X = (X_n)_{n \in \mathbb{Z}}$  be a stationary stochastic process. If*

$$\sup_{N \geq 1} \left( N \sum_{|n| \geq N} |\text{Cov}(X_0, X_n)| \right) < \infty, \quad (12)$$

and

$$\sum_{n \in \mathbb{Z}} \text{Cov}(X_0, X_n) = 0. \quad (13)$$

Then  $X$  is linearly rigid.



*Remark 3.1.* The condition (12) is a sufficient condition such that the spectral density  $\omega_X$  is a function in the Zygmund class  $\Lambda_*(1)$ , see below for definition. The condition (13) implies in particular that  $\omega_X$  vanishes at the point  $0 \in \mathbb{T}$ .

We shall apply a result of F. Móricz [12, Thm. 3] on absolutely convergent Fourier series and Zygmund class functions. Recall that a continuous 1-periodic function  $\varphi$  defined on  $\mathbb{R}$  is said to be in the *Zygmund class*  $\Lambda_*(1)$ , if there exists a constant  $C$  such that

$$|\varphi(x+h) - 2\varphi(x) + \varphi(x-h)| \leq Ch \quad (14)$$

for all  $x \in \mathbb{R}$  and for all  $h > 0$ .

**Theorem 3.2** (Móricz, [12]). *If  $\{c_n\}_{n \in \mathbb{Z}} \in \mathbb{C}$  is such that*

$$\sup_{N \geq 1} \left( N \sum_{|n| \geq N} |c_n| \right) < \infty, \quad (15)$$

*then the function  $\varphi(\theta) = \sum_{n \in \mathbb{Z}} c_n e^{i2\pi n\theta}$  is in the Zygmund class  $\Lambda_*(1)$ .*

*Proof of Theorem 3.1.* First, in view of (10), our assumption (13) implies

$$\omega_X(0) = 0.$$

Next, by Theorem 3.2, under the assumption (12), we have

$$\omega_X \in \Lambda_*(1).$$

Since all Fourier coefficients of  $\omega_X$  are real, we have

$$\omega_X(\theta) = \omega_X(-\theta).$$

Consequently, there exists  $C > 0$ , such that

$$\omega_X(\theta) = \frac{\omega_X(\theta) + \omega_X(-\theta)}{2} = \frac{\omega_X(\theta) + \omega_X(-\theta) - 2\omega_X(0)}{2} \leq C|\theta|,$$

whence

$$\int_{\mathbb{T}} \omega_X^{-1} dm = \infty,$$

and the stochastic process  $X = (X_n)_{n \in \mathbb{Z}}$  is linearly rigid by the Kolmogorov criterion.  $\square$

## 4 Applications to stationary determinantal point processes

In this section, we first give a sufficient condition for linear rigidity of stationary determinantal point processes on  $\mathbb{R}$  and then give an example of a very simple stationary, but not linearly rigid, determinantal point process on  $\mathbb{R}^2$ . We briefly recall the main definitions. Let  $B \subset \mathbb{R}^d$  be a bounded Borel subset. Let  $K_B : L^2(\mathbb{R}^d) \rightarrow L^2(\mathbb{R}^d)$  be the operator of convolution with the Fourier transform  $\widehat{\chi_B}$  of the indicator function  $\chi_B$ . In other words, the kernel of  $K_B$  is

$$K_B(x, y) = \widehat{\chi_B}(x - y). \quad (16)$$

In particular, if  $d = 1$  and  $B = (-1/2, 1/2)$ , then we find the well-known Dyson sine kernel

$$K_{\text{sine}}(x, y) = \frac{\sin(\pi(x - y))}{\pi(x - y)}.$$

Note that we always have  $K_B(x, x) = K_B(0, 0)$ .

Denote by  $\mathbb{P}_{K_B}$  the determinantal point process induced by  $K_B$ . For the background on the determinantal point processes, the reader is referred to [8], [9], [11], [13].

**Proposition 4.1.** *Let  $\mathbb{P}_{K_B}$  be the stationary determinantal point process on  $\mathbb{R}^d$  induced by the kernel  $K_B$  in (16). For any  $\lambda > 0$ , denote by  $N^{(\lambda)} = (N_n^{(\lambda)})_{n \in \mathbb{Z}^d}$  the stationary stochastic process associated to  $\mathbb{P}_{K_B}$  as in (2). Then*

$$\sum_{n \in \mathbb{Z}^d} |\text{Cov}(N_0^{(\lambda)}, N_n^{(\lambda)})| < \infty \quad (17)$$

and

$$\sum_{n \in \mathbb{Z}^d} \text{Cov}(N_0^{(\lambda)}, N_n^{(\lambda)}) = 0. \quad (18)$$

*Proof.* Fix a number  $\lambda > 0$ , for simplifying the notation, let us denote  $N_n^{(\lambda)}$  by  $N_n$ . Denote for any  $n \in \mathbb{Z}^d$ ,

$$Q_n = n\lambda + [-\lambda/2, \lambda/2]^d.$$

By definition of a determinantal point process, we have

$$\mathbb{E}(N_n) = \mathbb{E}(N_0) = \int_{Q_0} K_B(x, x) dx = \lambda^d K_B(0, 0).$$

If  $n \neq 0$ , we have

$$\begin{aligned}\mathbb{E}(N_0 N_n) &= \iint \chi_{Q_0}(x) \chi_{Q_n}(y) \begin{vmatrix} K_B(x, x) & K_B(x, y) \\ K_B(y, x) & K_B(y, y) \end{vmatrix} dx dy \\ &= \lambda^{2d} K_B(0, 0)^2 - \iint_{Q_0 \times Q_n} |K_B(x, y)|^2 dx dy,\end{aligned}$$

whence

$$\text{Cov}(N_0, N_n) = - \iint_{Q_0 \times Q_n} |K_B(x, y)|^2 dx dy. \quad (19)$$

We also have

$$\begin{aligned}\mathbb{E}(N_0^2) &= \mathbb{E} \left[ \sum_{x, y \in \mathcal{X}} \chi_{Q_0}(x) \chi_{Q_0}(y) \right] \\ &= \mathbb{E} \left[ \sum_{x \in \mathcal{X}} \chi_{Q_0}(x) \right] + \mathbb{E} \left[ \sum_{x, y \in \mathcal{X}, x \neq y} \chi_{Q_0}(x) \chi_{Q_0}(y) \right] \\ &= \int_{Q_0} K_B(x, x) dx + \iint \chi_{Q_0}(x) \chi_{Q_0}(y) \begin{vmatrix} K_B(x, x) & K_B(x, y) \\ K_B(y, x) & K_B(y, y) \end{vmatrix} dx dy \\ &= \lambda^d K_B(0, 0) + \lambda^{2d} K_B(0, 0)^2 - \iint_{Q_0 \times Q_0} |K_B(x, y)|^2 dx dy,\end{aligned}$$

whence

$$\text{Cov}(N_0, N_0) = \text{Var}(N_0) = \lambda^d K_B(0, 0) - \iint_{Q_0 \times Q_0} |K_B(x, y)|^2 dx dy. \quad (20)$$

Now recall that  $K_B$  is an orthogonal projection. Thus we have

$$K_B(0, 0) = K_B(x, x) = \int |K_B(x, y)|^2 dy = \sum_{n \in \mathbb{Z}^d} \int_{Q_n} |K_B(x, y)|^2 dy. \quad (21)$$

The identities (19), (20) and (21) imply that

$$\begin{aligned}\sum_{n \in \mathbb{Z}^d} \text{Cov}(N_0, N_n) &= \lambda^d K_B(0, 0) - \int_{Q_0} dx \sum_{n \in \mathbb{Z}^d} \int_{Q_n} |K_B(x, y)|^2 dy \\ &= \lambda^d K_B(0, 0) - \lambda^d K_B(0, 0) = 0.\end{aligned}$$

Moreover, the above series converge absolutely. Proposition 4.1 is completely proved.  $\square$

*Remark 4.1.* By Lemma 2.3 and Proposition 4.1, we see that for any stationary determinantal point process induced by a projection, the spectral density of the associated stochastic process  $N^{(\lambda)}$  always vanishes at 0.

## 4.1 Stationary determinantal point processes on $\mathbb{R}$

**Theorem 4.2.** Assume that  $B \subset \mathbb{R}$  satisfies

$$\sup_{R>0} \left( R \int_{|\xi| \geq R} |\widehat{\chi}_B(\xi)|^2 d\xi \right) < \infty. \quad (22)$$

Then the stationary determinantal point process  $\mathbb{P}_{K_B}$  is linearly rigid.

*Proof.* By definition of linear rigidity, we need to show that for any  $\lambda > 0$ , the stochastic process  $N^{(\lambda)} = (N_n^{(\lambda)})_{n \in \mathbb{Z}}$  is linearly rigid. As in the proof of Proposition 4.1, we denote  $N_n^{(\lambda)}$  by  $N_n$ . By Theorem 3.1, it suffices to show that

$$\sup_{N \geq 1} \left( N \sum_{|n| \geq N} |\text{Cov}(N_0, N_n)| \right) < \infty, \quad (23)$$

and

$$\sum_{n \in \mathbb{Z}} \text{Cov}(N_0, N_n) = 0. \quad (24)$$

By Proposition 4.1, the identity (24) holds in general case. It remains to prove (23). By (19), we have

$$\begin{aligned} \sup_{N \geq 1} \left( N \sum_{|n| \geq N} |\text{Cov}(N_0, N_n)| \right) &= \sup_{N \geq 1} N \iint_{\bigcup_{|n| \geq N} Q_n} |\widehat{\chi}_B(x - y)|^2 dx dy \\ &\leq \sup_{N \geq 1} \lambda N \int_{|\xi| \geq (N-1)\lambda} |\widehat{\chi}_B(\xi)|^2 d\xi < \infty \end{aligned}$$

where in the last inequality, we used our assumption (22). Theorem 4.2 is proved completely.  $\square$

*Remark 4.2.* When  $B$  is a finite union of finite intervals on the real line, the rigidity of the stationary determinantal point process  $\mathbb{P}_{K_B}$  is due to Ghosh [5].

## 4.2 Tensor product of sine kernels

In higher dimension, the situation becomes quite different. Let

$$S = I \times I = (-1/2, 1/2) \times (-1/2, 1/2) \subset \mathbb{R}^2.$$

Then the associate kernel  $K_S$  has a tensor form:  $K_S = K_{\text{sine}} \otimes K_{\text{sine}}$ , that is, for  $x = (x_1, x_2)$  and  $y = (y_1, y_2)$  in  $\mathbb{R}^2$ , we have

$$K_S(x, y) = K_{\text{sine}}(x_1, y_1)K_{\text{sine}}(x_2, y_2) = \frac{\sin(\pi(x_1 - y_1))}{\pi(x_1 - y_1)} \frac{\sin(\pi(x_2 - y_2))}{\pi(x_2 - y_2)}.$$

**Proposition 4.3.** *The determinantal point process  $\mathbb{P}_{K_S}$  is not linearly rigid. More precisely, let  $N^{(1)} = (N_n^{(1)})_{n \in \mathbb{Z}^2}$  be the stationary stochastic process given as in Definition 1.2, then*

$$N_0^{(1)} \notin \check{H}_0(N^{(1)}).$$

To prove the above result, we need to introduce some extra notation. First, we define the multiple Zygmund class  $\Lambda_*$  as follows. A continuous function  $\varphi(x, y)$  periodic in each variable with period 1 is said to be in the multiple Zygmund class  $\Lambda_*(1, 1)$  if for the double difference difference operator  $\Delta_{2,2}$  of second order in each variable, applied to  $\varphi$ , there exists a constant  $C > 0$ , such that for all  $x = (x_1, x_2) \in (-1/2, 1/2) \times (-1/2, 1/2)$  and  $h_1, h_2 > 0$ , we have

$$|\Delta_{2,2}\varphi(x_1, x_2; h_1, h_2)| \leq Ch_1h_2, \quad (25)$$

where

$$\begin{aligned} \Delta_{2,2}\varphi(x_1, x_2; h_1, h_2) &:= \varphi(x_1 + h_1, x_2 + h_2) + \varphi(x_1 - h_1, x_2 + h_2) \\ &\quad + \varphi(x_1 + h_1, x_2 - h_2) + \varphi(x_1 - h_1, x_2 - h_2) - 2\varphi(x_1 + h_1, x_2) \\ &\quad - 2\varphi(x_1 - h_1, x_2) - 2\varphi(x_1, x_2 + h_2) - 2\varphi(x_1, x_2 - h_2) + 4\varphi(x_1, x_2). \end{aligned}$$

The following result is due to Fülöp and Móricz [4, Thm 2.1 and Rem. 2.3]

**Theorem 4.4** (Fülöp-Móricz). *If  $\{c_{jk}\}_{j,k \in \mathbb{Z}} \in \mathbb{C}$  is such that*

$$\sup_{N \geq 1, M \geq 1} \left( MN \sum_{|j| \geq N, |k| \geq M} |c_{jk}| \right) < \infty, \quad (26)$$

*then the function*

$$\varphi(\theta_1, \theta_2) = \sum_{j,k \in \mathbb{Z}} c_{jk} e^{i2\pi(j\theta_1 + k\theta_2)}$$

*is in the Zygmund class  $\Lambda_*(1, 1)$ .*

Let us turn to the study of the density function  $\omega_{N^{(1)}}$ .

**Lemma 4.5.** *There exists  $c > 0$ , such that for any  $\theta_1, \theta_2 \in (-1/2, 1/2)$ , we have*

$$\omega_{N^{(1)}}(\theta_1, \theta_2) \geq c(|\theta_1| + |\theta_2|).$$

*Proof.* To make notation lighter, in this proof we simply write  $\omega$  for  $\omega_{N^{(1)}}$ .

Denote  $S_n = S \times (n + S)$  where  $n + S := (-1/2 + n_1, 1/2 + n_1) \times (-1/2 + n_2, 1/2 + n_2)$ . By the same argument as in the proof of Theorem 4.2, we obtain that for any  $n = (n_1, n_2) \in \mathbb{Z}^2 \setminus 0$ ,

$$\widehat{\omega}(n) = - \int_{S_n} |K_S(x, y)|^2 dx dy,$$

and

$$\widehat{\omega}(0) = K_S(0, 0) - \int_{S_0} |K_S(x, y)|^2 dx dy.$$

The following properties can be easily checked.

- $\sum_{n \in \mathbb{Z}^2} \widehat{\omega}(n) = 0$ .
- $\widehat{\omega}(\varepsilon_1 n_1, \varepsilon_2 n_2) = \widehat{\omega}(n_1, n_2)$ , where  $\varepsilon_1, \varepsilon_2 \in \{\pm 1\}$ .
- there exist  $c, C > 0$ , such that

$$\frac{c}{(1 + n_1^2)(1 + n_2^2)} \leq |\widehat{\omega}(n_1, n_2)| \leq \frac{C}{(1 + n_1^2)(1 + n_2^2)}.$$

For instance,  $\sum_{n \in \mathbb{Z}^2} \widehat{\omega}(n) = 0$  follows from Proposition 4.1. These properties combined with Theorem 4.4 yield that

- $\omega(0, 0) = 0$ .
- $\omega(\varepsilon_1 \theta_1, \varepsilon_2 \theta_2) = \omega(\theta_1, \theta_2)$  for any  $\varepsilon_1, \varepsilon_2 \in \{\pm 1\}$  and  $\theta_1, \theta_2 \in (-1/2, 1/2)$ .
- the function  $\omega(\theta_1, \theta_2)$  is in the multiple Zygmund class  $\Lambda_*(1, 1)$ .

Hence there exists  $C > 0$ , such that

$$|\omega(\theta_1, \theta_2) - \omega(\theta_1, 0) - \omega(0, \theta_2)| \leq C|\theta_1 \theta_2|. \quad (27)$$

**Lemma 4.6.** *There exists  $c > 0$ , such that*

$$\omega(\theta_1, 0) \geq c|\theta_1| \text{ and } \omega(0, \theta_2) \geq c|\theta_2|. \quad (28)$$

Let us postpone the proof of Lemma 4.6 and proceed to the proof of Lemma 4.5. The inequalities (27) and (28) imply that

$$\omega(\theta_1, \theta_2) \geq c(|\theta_1| + |\theta_2|) - C|\theta_1\theta_2|.$$

To prove the lower bound of type as in the lemma, it suffices to prove it when  $|\theta_1|$  and  $|\theta_2|$  are small enough, for instance,  $2C|\theta_1| \leq c$ , then we have

$$\omega(\theta_1, \theta_2) \geq \frac{c}{2}(|\theta_1| + |\theta_2|).$$

□

Now let us turn to the

*Proof of Lemma 4.6.* By symmetry, it suffices to prove that there exists  $c > 0$ , such that  $\omega(\theta_1, 0) \geq |\theta_1|$ . To this end, let us denote  $\omega_1(\theta_1) := \omega(\theta_1, 0)$ . Then  $\omega_1(0) = 0$  and there exists  $c > 0$  such that if  $k \neq 0$ , then

$$\widehat{\omega}_1(k) < 0 \text{ and } |\widehat{\omega}_1(k)| \geq c/(1 + k^2),$$

Indeed, we have

$$\omega_1(\theta_1) = \sum_{k \in \mathbb{Z}} \sum_{n_2 \in \mathbb{Z}} \widehat{\omega}(k, n_2) e^{i2\pi k \theta_1},$$

if  $k \neq 0$ , then  $\widehat{\omega}(k, n_2) < 0$  and hence

$$|\widehat{\omega}_1(k)| = \sum_{n_2 \in \mathbb{Z}} |\widehat{\omega}(k, n_2)| \geq \sum_{n_2 \in \mathbb{Z}} \frac{c}{(1 + n_2^2)(1 + k^2)} \geq \frac{c'}{1 + k^2}.$$

Note also that  $\omega_1(0) = \omega(0, 0) = 0$ , hence

$$\sum_{k \in \mathbb{Z}} \widehat{\omega}_1(k) = 0.$$

It follows that

$$\begin{aligned}
\omega_1(\theta_1) &= \sum_{k \in \mathbb{Z}} \widehat{\omega}_1(k) e^{i2\pi k \theta_1} = \sum_{k \in \mathbb{Z}} \widehat{\omega}_1(k) \left( \frac{e^{i2\pi k \theta_1} + e^{-i2\pi k \theta_1}}{2} - 1 \right) \\
&= \sum_{k \in \mathbb{Z}, k \neq 0} -\widehat{\omega}_1(k) (1 - \cos(2\pi k \theta_1)) = \sum_{k \in \mathbb{Z}, k \neq 0} |\widehat{\omega}_1(k)| (1 - \cos(2\pi k \theta_1)) \\
&\geq c'' \sum_{j=1}^{\infty} \frac{1}{(2j-1)^2} (1 - \cos(2\pi(2j-1)\theta_1)).
\end{aligned}$$

Combining with the classical formulae

$$\begin{aligned}
\sum_{j=1}^{\infty} \frac{1}{(2j-1)^2} &= \frac{\pi^2}{8}, \\
|\alpha| &= \frac{1}{4} - \frac{2}{\pi^2} \sum_{j=1}^{\infty} \frac{\cos(2(2j-1)\pi\alpha)}{(2j-1)^2}, \text{ for } \alpha \in (-1/2, 1/2);
\end{aligned}$$

we obtain that

$$\omega_1(\theta_1) \geq c'' \frac{\pi^2}{2} |\theta_1|.$$

□

*Proof of Proposition 4.3.* By Lemma 2.1, it suffices to show that

$$\int_{\mathbb{T}^2} \omega_{N(1)}^{-1} dm < \infty. \quad (29)$$

By Lemma 4.5, the inequality (29) follows from the following elementary inequality

$$\int_{|\theta_1| < 1/2, |\theta_2| < 1/2} \frac{1}{|\theta_1| + |\theta_2|} d\theta_1 d\theta_2 < \infty.$$

□



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